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# Optimal modeling of nonlinear systems: Method of variable injections 

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#### Abstract

Our work addresses a development and justification of the new approach to the modeling of nonlinear systems. Let $\mathcal{F}$ be an unknown input-output map of the system with a random input and output $\mathbf{Y}$ and $\mathbf{X}$, respectively. It is assumed that $\mathbf{Y}$ and $\mathbf{X}$ are available and covariance matrices formed from $\mathbf{Y}$ and $\mathbf{X}$ are known. We determine a model of $\mathcal{F}$ so that an associated error is minimized. To this end, the model $\mathcal{T}_{p}$ is constructed as a sum of $p+1$ particular parts, in the form $\mathcal{T}_{p}(\mathbf{y})=\sum_{j=0}^{p} G_{j} H_{j} Q_{j}\left(\mathbf{V}_{j}\right)$ where $G_{j}$ and $H_{j}$, for $j=0, \ldots, p$, are matrices to be determined, and $\mathbf{V}_{j}$, for $j=1, \ldots, p$, is a special random vector called the injection. We denote $\mathbf{V}_{0}=\mathbf{y}$. Further, $Q_{j}$ is a special transform aimed to facilitate the numerical realization of model $\mathcal{T}_{p}$. It is determined in the way allowing us to optimally determine $G_{j}$ and $H_{j}$ as a solution of $p+1$ separate error minimization problems which are simpler than the original minimization problem. The empirical determination of injections $\mathbf{V}_{1}, \ldots, \mathbf{V}_{p}$ is considered. The proposed method has several degrees of freedom to diminish the associated error. They are 'degree' $p$ of $\mathcal{T}_{p}$, choice of matrices $G_{0}, H_{0}, \ldots$, $G_{p}, H_{p}$, dimensions of matrices $G_{0}, H_{0}, \ldots, G_{p}, H_{p}$ and injections $\mathbf{V}_{1}, \ldots, \mathbf{V}_{p}$, respectively. In particular, it is shown that a variation of the injections in their dimensionality and special forms allow us to increase accuracy of the proposed model $\mathcal{T}_{p}$. The proposed approach differs from known techniques by its ingredients mentioned above. Four numerical examples are provided. At the end, the open problem is formulated.


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## 1. Introduction

We propose and justify a new method for the optimal modeling of systems transforming random signals.

### 1.1. Motivation

The problem of a constructive representation of nonlinear systems has been a topic of profound research for a long time. A number of fundamental papers have appeared which established significant advances in this research area. Some relevant references can be found, in particular, in $[31,19,67$, $69,68,34,24,15,22,21,72,6,20,63,50,42]$.

The known related results mainly concern proving the existence and uniqueness of mathematical models used for an approximate representation of a system under consideration, and for justifying the bounds of errors arising from the approximation methods. The assumptions are that inputs and outputs are deterministic and can be represented in an analytical form, that is, by equations. At the same time, in many applications, the inputs and outputs are stochastic and cannot be described by equations. Nevertheless, it is possible to represent these sets in terms of their numerical characteristics, such as the expectation and covariance matrices. Typical examples are engineering $[66,44,23,8]$, statistics $[15,38,16,74]$, stochastic signal processing [61, 36, 92, 88, 9], and image processing [17, 60]; in the latter case, a digitized image, presented by a matrix, is often interpreted as the sample of a stochastic signal.

While the theory of a system representation with any given accuracy is well elaborated (see, e.g., [31, 19, 67, 69, 68, 34, 24]), the theory of optimal constrained and constructive system representation is still not so well developed, although this is an area of intensive research (see, e.g., [22, 21]). Despite increasing demands from applications [72, 6, 20, 63, 92, $88,9,38,16,74,66,44,23,8,60,22,21,28,79,51,3,32,29,5,62,65]$ this subject is hardly tractable because of intrinsic difficulties in optimal approximation techniques, especially when the approximating model should have a specific structure implied by the underlying problem.

We wish to extend the known results in this area to the case when the inputs and outputs of the system are stochastic, and the approximating model we search is constructive. The latter means the model can numerically be realized and, therefore, is applicable to problems in applications.

### 1.2. Differences from known techniques

For the case when the dimensionality reductions of signals are not required, the following results related to the problem under consideration are as follows.

The well-known Levinson-Schur algorithms (see, e.g. [78, 77]) are fast but they are only applicable to Toeplitz matrices. The covariance matrices we consider are not Toeplitz. It is possible to replace the covariance matrix with a Toeplitz matrix that is closest to a given covariance matrix [26], and then apply the Levinson-Schur algorithms to that. But this replacement will, of course, increase the associated error.

The alternating direction method of multipliers (ADMM) is well suited to the distributed convex optimization [13] (see pp. 2, 33-37) but the problem we consider is not convex. For non-convex problems, ADMM can converge to non-optimal points - see pp. 73-77 in [13]. In [57], the proposed scenario requires an use of block diagonal matrices and the existence of non-singular matrices. In our solution, we do not limit ourselves with those restrictions. In [14], a fast low-rank modifications of the thin SVD is studied, but the SVD of the covariance matrices we consider is not thin. Methods of a fast computation of the pseudo-inverse matrix in [18, 7, 55] are approximate and iterative and, therefore, the associated computational load increases with the increase in the number of iterations. The proposed method avoids iterations for the determination of the optimal system model.

We also note that the problems we consider are different from those studied in [86, 54, 41].

The optimal dimensionality reduction of a random signal is achieved by the Karhunen-Loève transform (KLT) which is studied in a number of works. We cite just a few of them in [70, 35, 36]. The KLT is also known as Principal Component Analysis [37]. Recall, the KLT is represented by a matrix of fixed rank that minimizes the error associated with the optimal compression and estimation of the source random vector $s$. We consider a case where the system model consists of $p+1$ blocks transforming input signal $\mathbf{y}$. If the number of blocks $p+1$ is greater than one then the compression of data arrays cannot be done by the KLT. Therefore, to find a model of the system, for $p>1$, a new approach should be used. In particular, for the case of two system blocks, this problem has been solved in [83].

For an arbitrary number of the system blocks, in $[27,25,75,93,71,4$, $48,45,56,89,91]$, approximate iterative solutions of this problem are considered. They are mainly based on implementations of the block coordinate descent method (BCDM) and its modifications [84, 10]. Unfortunately, con-
vergence of the BCDM to a global minimum is derived under some heavy restrictions $[84,10]$ which are difficult to implement and often cannot be satisfied in practice.

The proposed method is based on an approach which is different from those mentioned above. This allows us avoid associated difficulties.

### 1.3. Short description of the method

Let $\Omega$ be a set of outcomes in probability space $(\Omega, \Sigma, \mu)$ for which $\Sigma$ is a $\sigma$-field of measurable subsets of $\Omega$ and $\mu: \Sigma \rightarrow[0,1]$ is an associated probability measure. Let $\mathbf{x} \in L^{2}\left(\Omega, \mathbf{R}^{m}\right)$ and $\mathbf{y} \in L^{2}\left(\Omega, \mathbf{R}^{n}\right)$ be random input and output, respectively, such that $\mathbf{x}=\mathcal{F}(\mathbf{y})$ where $\mathcal{F}: L^{2}\left(\Omega, \mathbf{R}^{n}\right) \rightarrow$ $L^{2}\left(\Omega, \mathbf{R}^{m}\right)$ is an unknown non-linear input-output map. It is assumed that the only available information on $\mathcal{F}$ is given by certain covariance matrices formed from $\mathbf{x}$ and $\mathbf{y}$. This is a typical assumption used in the applications such as those considered, e.g., in $[61,36,92,88,9,15,38,16,74,72,6,5$, $62,65]$. Here, we adopt that assumption. It is also assumed that $\mathbf{x}$ and $\mathbf{y}$ are available.

We develop a new approach to the optimal constructive representation of the nonlinear system $\mathcal{F}$ subject to a specialized criterion associated with the dimensionality reduction of the random input. The latter constraint follows from the requirements in applications such as those considered in $[15,38,16,74,44,3,32]$. In particular, a dimensionality reduction of random signals is used to optimize the cost of signal transmission.

In Section 2.3 that follows, we consider a problem that concerns finding the optimal system model that depends on $2 p+2$ unknown matrices $G_{j}$ and $H_{j}$, for $j=0, \ldots, p$ and a given non-negative integer $p$. The difficulty is that $2 p+2$ unknown matrices should be determined from a minimization of the single cost function represented in (2.6). The solution is provided in Section 3 and is given by the optimal approximating operator minimizing the associated error.

Numerical examples that illustrate the advantages of the proposed method are provided. Computations were performed on a desktop computer with an $\operatorname{Intel}(\mathrm{R})$ Core(TM) i9-10900F CPU ( 2.80 GHz ) and 32.00 RAM using Matlab R2021a. The related code is available in GitHub. The link is https://github.com/jusotoTEC/variableInjections.

### 1.4. Novelty

The proposed system model $\mathcal{T}_{p}$ is represented by a sum of $p+1$ terms where $j$ th term, for $j=0, \ldots, p$, is a model of $j$ th block of the system (see (2.5) in Section 2.2). Each term of model $\mathcal{T}_{p}$ contains a special transformation $Q_{j}$, for $j=0, \ldots, p$, which leads to a faster numerical realization of the proposed method. This is because the transformations provide a representation of the problem under consideration in (2.6) by a set of simpler problem each of which depends on the single pair $G_{j}, H_{j}$, for $j=0, \ldots, p$. It allows us to avoid numerical difficulties associated with computation of large matrices and determine $G_{j}, H_{j}$, for $j=0, \ldots, p$ from the problem which is simpler than that in (2.6). Details are given in Section 3.2 that follows.

The model $\mathcal{T}_{p}$ has several degrees of freedom to minimize the associated error. They are: 'degree' $p$ of $\mathcal{T}_{p}$, choices of matrices $G_{0}, H_{0}, \ldots, G_{p}, H_{p}$ (Theorems 3, 4, 5) and injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ (Section 3.3, Theorem 5, Remarks 2,3 and 6 ), and dimensions $r_{0}, \ldots, r_{p}$ (Theorems 3). In particular, an influence of variations of the injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ on the decrease of the associated error is shown in Theorem 4, and Remarks 2, 3 and 6: an increase in the dimensions and of the injections and their choice improve the accuracy of the system model.

## 2. The proposed approach

### 2.1. Some special notation

Let us write $\mathbf{x}=\left[\mathbf{x}_{(1)}, \ldots, \mathbf{x}_{(m)}\right]^{T}$ and $\mathbf{y}=\left[\mathbf{y}_{(1)}, \ldots, \mathbf{y}_{(n)}\right]^{T}$ where $\mathbf{x}_{(i)}, \mathbf{y}_{(j)} \in$ $L^{2}(\Omega, \mathbf{R})$, for $i=1, \ldots, m$ and $j=1, \ldots, n$, and $\mathbf{x}(\omega) \in R^{m}$ and $\mathbf{y}(\omega) \in R^{n}$ for all $\omega \in \Omega$.

Each matrix $A \in R^{m \times n}$ defines a bounded linear transformation $\mathcal{A}: L^{2}\left(\Omega, \mathbf{R}^{n}\right) \rightarrow L^{2}\left(\Omega, \mathbf{R}^{m}\right)$. It is customary to write $A$ rather then $\mathcal{A}$ since $[\mathcal{A}(\mathbf{x})](\omega)=A[\mathbf{x}(\omega)]$, for each $\omega \in \Omega$.

Let us also denote

$$
\begin{equation*}
\|\mathbf{x}\|_{\Omega}^{2}=\int_{\Omega} \sum_{j=1}^{m}\left[\mathbf{x}_{j}(\omega)\right]^{2} d \mu(\omega)<\infty \tag{2.1}
\end{equation*}
$$

The covariance matrix formed from $\mathbf{x}$ and $\mathbf{y}$ is denoted by

$$
\begin{equation*}
E_{x y}=\left\{\int_{\Omega} \mathbf{x}_{(i)}(\omega) \mathbf{y}_{(j)}(\omega) d \mu(\omega)\right\}_{i, j=1}^{m, n} \tag{2.2}
\end{equation*}
$$

The Moore-Penrose pseudo-inverse [11] of matrix $M$ is denoted by $M^{\dagger}$.

### 2.2. Generic structure of the system model

Let $\mathbf{y}$ be an input signal and $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ be random vectors such that $\mathbf{v}_{j} \in L^{2}\left(\Omega, \mathbf{R}^{q_{j}}\right)$, for $j=1, \ldots, p$. We write $\mathbf{y}=\mathbf{v}_{0}$ and $q_{0}=n$. Each $\mathbf{v}_{j}$, for $j=1, \ldots, p$, is defined by a nonlinear transformation $\varphi_{j}$ of $\mathbf{y}$, i.e., $\mathbf{v}_{j}=$ $\varphi_{j}(\mathbf{y})$. The choice of $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ is considered in Section 3.3. To facilitate the numerical implementation of the system model introduced below, each vector $\mathbf{v}_{j}$, for $j=1, \ldots, p$, is transformed to vector $\mathbf{z}_{j} \in L^{2}\left(\Omega, \mathbf{R}^{q_{j}}\right)$ by transformation $Q_{j}$ so that

$$
\begin{equation*}
\mathbf{z}_{j}=Q_{j}\left(\mathbf{v}_{j}, Z_{j-1}\right) \tag{2.3}
\end{equation*}
$$

where $Z_{j-1}=\left\{\mathbf{z}_{0}, \ldots, \mathbf{z}_{j-1}\right\}$. The choice of $Q_{j}$ is considered in Section 3.1. We call $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ the injections. This is because $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ contribute to the decrease of the associated error as shown in Section 3.4 below. At the same time, in particular, $\mathbf{v}_{j}$ can be a noisy version of $\mathbf{y}$, e.g., $\mathbf{v}_{j}=\mathbf{y}+\mathbf{n}_{j}$ where $\mathbf{n}_{j}$ is a random noise. In this case, system filters signals $\mathbf{v}_{0}, \ldots, \mathbf{v}_{p}$.

Further, for $i=0,1, \ldots, p$, let $G_{i} \in R^{m \times r_{i}}, H_{i} \in R^{r_{i} \times q_{i}}$ where $r_{i}$ is given, $0<r_{i}<r$ and

$$
\begin{equation*}
r=r_{0}+\ldots+r_{p} \tag{2.4}
\end{equation*}
$$

Here, $r$ is a positive integer such that $r \leq \min \{m, n\}$. It is convenient to set $Q_{0}=I$ and $\mathbf{z}_{0}=\mathbf{v}_{0}=\mathbf{y}$.

For a given reduction ratio

$$
c=r / \min \{m, n\}
$$

we consider a system model given by operator $\mathcal{T}_{p}$ : $L^{2}\left(\Omega, \mathbf{R}^{q_{0}}\right) \times \ldots \times$ $L^{2}\left(\Omega, \mathbf{R}^{q_{p}}\right) \rightarrow L^{2}\left(\Omega, \mathbf{R}^{m}\right)$ such that

$$
\begin{equation*}
\mathcal{T}_{p}\left(\mathbf{v}_{0}, \ldots \mathbf{v}_{p}\right)=G_{0} H_{00}+\ldots+G_{p} H_{p} \mathbf{z}_{p} \tag{2.5}
\end{equation*}
$$

where $G_{j}: L^{2}\left(\Omega, \mathbf{R}^{r_{j}}\right) \rightarrow L^{2}\left(\Omega, \mathbf{R}^{m}\right)$ and $H_{j}: L^{2}\left(\Omega, \mathbf{R}^{q_{j}}\right) \rightarrow L^{2}\left(\Omega, \mathbf{R}^{r_{j}}\right)$, for $j=0,1, \ldots, p$, are linear operators (i.e. $G_{j}$ and $H_{j}$ are represented by $m \times r_{j}$ and $r_{j} \times q_{j}$ matrices, respectively. Recall, we use the same symbol to define a matrix and the associated liner operator).

Importantly, operators $H_{0}, \ldots, H_{p}$ imply the dimensionality reduction of vectors $\mathbf{v}_{0}, \ldots, \mathbf{v}_{p}$ because $H_{i i} \in L^{2}\left(\Omega, \mathbf{R}^{r_{i}}\right)$ where $0<r_{i}<r \leq$ $\min \{m, n\}$, for $i=0, \ldots, p$.

We call $p$ the degree of model $\mathcal{T}_{p}$. It is shown below that $\mathcal{T}_{p}$ approximates the system of interest $\mathcal{F}: L^{2}\left(\Omega, \mathbf{R}^{n}\right) \rightarrow L^{2}\left(\Omega, \mathbf{R}^{m}\right)$ with the accuracy represented by Theorems 8-6 below.

### 2.3. Statement of the problem

Let $\mathcal{F}: L^{2}\left(\Omega, \mathbf{R}^{n}\right) \rightarrow L^{2}\left(\Omega, \mathbf{R}^{m}\right)$ be a continuous map. We consider the problem as follows: Given $\mathbf{x}, \mathbf{y}, \mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ and $r_{0}, \ldots, r_{p}$, find matrices $G_{0}, H_{0}, \ldots, G_{p}, H_{p}$ and transformations $Q_{0}, \ldots, Q_{p}$ that solve

$$
\begin{equation*}
\min _{G_{0}, H_{0}, \ldots, G_{p}, H_{p}}\left\|\mathcal{F}(\mathbf{y})-\sum_{j=0}^{p} G_{j} H_{j} \mathbf{z}_{j}\right\|_{\Omega}^{2} \tag{2.6}
\end{equation*}
$$

subject to

$$
\begin{equation*}
G_{j} \in R^{m \times r_{j}} \text { and } H_{j} \in R^{r_{j} \times q_{j}} \tag{2.7}
\end{equation*}
$$

and

$$
\begin{equation*}
E_{z_{i} z_{j}}=\mathbf{O}, \quad \text { for } i \neq j \tag{2.8}
\end{equation*}
$$

where $i, j=0, \ldots, p$ and $\mathbf{O}$ denotes the zero matrix (and the zero vector). Here, $\mathbf{z}_{j}$, for $j=0, \ldots, p$, is represented by (2.3). We note that, for $F_{j}=$ $G_{j} H_{j}$, the constraint in (2.7) can also be written as

$$
\operatorname{rank} F_{j}=r_{j} .
$$

It will be shown in Section 3 below that the solution of problem (2.6) - (2.8) is determined under a special condition imposed on vectors $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$. The condition is provided by Definition 3.

An assumption used in the methods of transformations of random signals is that the covariance matrices $E_{x y}$ and $E_{y y}$ are known (see, for example, $[36,15,40,65,82])$. Here, we adopt this assumption. As mentioned, in particular, in [46, 47, 2, 85, 39, 90], a priori knowledge of the covariances can come either from specific data models, or, after sample estimation during a training phase. Examples 1 and 3 in Sections 3.3 and 3.4 below illustrate this observation.

## 3. Solution of problem (2.6)-(2.8)

Definition 1. Random vectors $\mathbf{z}_{0}, \ldots, \mathbf{z}_{p}$ are called pairwise uncorrelated if the condition in (2.8) holds for any pair of vectors $\mathbf{z}_{i}$ and $\mathbf{z}_{j}$, for $i \neq j$, where $i, j=0, \ldots, p$. Two vectors $\mathbf{z}_{i}$ and $\mathbf{z}_{j}$ belonging to the set of the pairwise uncorrelated vectors are called uncorrelated.

For $j=0, \ldots, p$, let $\mathcal{N}\left(M^{(j)}\right)$ be a null space of matrix $M^{(j)} \in R^{q_{j} \times q_{j}}$. Definition 2. Random vectors $\mathbf{v}_{0}, \ldots, \mathbf{v}_{p}$ are called jointly independent if

$$
M^{(0)} \mathbf{v}_{0}(\omega)+\ldots+M^{(p)} \mathbf{v}_{p}(\omega)=\mathbf{O}
$$

almost everywhere in $\Omega$, only if $\mathbf{v}_{j}(\omega) \in \mathcal{N}\left(M^{(j)}\right)$, for $j=0, \ldots, p$.
Definition 3. Random vector $\mathbf{v}_{j}$, for $j=0, \ldots, p$, is called the well-defined injection if

$$
\begin{equation*}
\Gamma_{z_{j}}=E_{x z_{j}} E_{z_{j} z_{j}}^{\dagger} E_{z_{j} x} \neq \mathbf{O} \tag{3.1}
\end{equation*}
$$

where $\mathbf{z}_{\mathbf{j}}$ is defined by (2.3). Otherwise, injection $\mathbf{v}_{j}$ is called ill-defined.
An explanation for introducing Definition 3 is provided by Remark 1 at the end of Section 3.2 below.

### 3.1. Determination of pairwise uncorrelated vectors

Theorem 1. Let random vectors $\mathbf{v}_{0}, \ldots, \mathbf{v}_{p}$ be jointly independent. Then they are transformed to the pairwise uncorrelated vectors $\mathbf{z}_{0}, \ldots, \mathbf{z}_{p}$ by transformations $Q_{0}, \ldots, Q_{p}$ as follows:

$$
\begin{align*}
& \mathbf{z}_{0}=Q_{0}\left(\mathbf{v}_{0}\right)=\mathbf{v}_{0} \quad \text { and, for } j=1, \ldots, p,  \tag{3.2}\\
& \mathbf{z}_{j}=Q_{j}\left(\mathbf{v}_{j}, Z_{j-1}\right)=\mathbf{v}_{j}-\sum_{k=0}^{j-1} E_{v_{j} z_{k}} E_{z_{k} z_{k}}^{\dagger} \mathbf{z}_{k} \tag{3.3}
\end{align*}
$$

Proof. Suppose that the condition in (2.8) holds for $\mathbf{z}_{0}, \ldots, \mathbf{z}_{i-1}$. Then, for $\ell=0, \ldots, i-1$,

$$
\begin{align*}
E_{z_{i} z_{\ell}} & =E\left[\left(\mathbf{v}_{i}-\sum_{l=0}^{i-1} E_{v_{i} z_{l}} E_{z_{l} z_{l}}^{\dagger} \mathbf{z}_{l}\right) \mathbf{z}_{\ell}^{T}\right] \\
& =E_{v_{i} z_{\ell}}-\sum_{l=0}^{i-1} E_{v_{i} z_{l}} E_{z_{l} z_{l}}^{\dagger} E_{z_{l} z_{\ell}} \\
& =E_{v_{i} z_{\ell}}-E_{v_{i} z_{\ell}} E_{z_{\ell} z_{\ell}}^{\dagger} E_{z_{\ell} z_{\ell}}=\mathbf{O} . \tag{3.4}
\end{align*}
$$

The latter is true because by Lemma 1 in [80],

$$
E_{v_{i} z_{\ell}} E_{z_{\ell} z_{\ell}}^{\dagger} E_{z_{\ell} z_{\ell}}=E_{v_{i} z_{\ell}} .
$$

Thus, by induction, (2.8) holds for any $i=0, \ldots, p$.

### 3.2. Determination of matrices $G_{0}, H_{0}, \ldots, G_{p}, H_{p}$ that solve problem (2.6)

First, recall the definition of a truncated SVD. Let the SVD of matrix $A \in R^{m \times n}$ be given by

$$
A=U_{A} \Sigma_{A} V_{A}^{T}
$$

where $U_{A}=\left[u_{1} u_{2} \ldots u_{m}\right] \in R^{m \times m}, V_{A}=\left[v_{1} v_{2} \ldots v_{n}\right] \in R^{n \times n}$ are unitary matrices, and $\Sigma_{A}=\left(\sigma_{1}(A), \ldots, \sigma_{\min (m, n)}(A)\right) \in R^{m \times n}$ is a generalized diagonal matrix, with the singular values $\sigma_{1}(A) \geq \sigma_{2}(A) \geq \ldots \geq 0$ on the main diagonal. For $k<m, j<n$ and $\ell<\min (m, n)$, we denote $U_{A, k}=\left[\begin{array}{llll}u_{1} & u_{2} & \ldots & u_{k}\end{array}\right], V_{A, j}=\left[\begin{array}{llll}v_{1} & v_{2} & \ldots & v_{j}\end{array}\right], \Sigma_{A, \ell}=\left(\sigma_{1}(A), \ldots, \sigma_{\ell}(A)\right)$, and write

$$
\Pi_{A, L}=\sum_{k=1}^{\operatorname{rank}(A)} u_{k} u_{k}^{T} \text { and } \Pi_{A, R}=\sum_{j=1}^{\operatorname{rank}(A)} v_{j} v_{j}^{T} .
$$

For $r=1, \ldots, \operatorname{rank}(A)$,

$$
\begin{equation*}
[A]_{r}=\sum_{i=1}^{r} \sigma_{i}(A) u_{i} v_{i}^{T} \in R^{m \times n} \tag{3.5}
\end{equation*}
$$

is the truncated SVD of $A$. For $r \geq \operatorname{rank}(A)$ we write $[A]_{r}=A(=$ $\left.A_{\text {rank (A) }}\right)$.

Theorem 2. Let $\mathbf{v}_{0}, \ldots, \mathbf{v}_{p}$ be well-defined injections and vectors $\mathbf{z}_{0}, \ldots, \mathbf{z}_{p}$ be pairwise uncorrelated. Then the minimal Frobenius norm solution to the problem in (2.6) is given, for $j=0, \ldots, p$, by

$$
\begin{equation*}
G_{j}=U_{\Gamma_{z_{j}}, r_{j}} \text { and } H_{j}=U_{\Gamma_{z_{j}}, r_{j}}^{T} E_{x z_{j}} E_{z_{j} z_{j}}^{\dagger} \tag{3.6}
\end{equation*}
$$

Proof. For $j=0, \ldots, p$, let $S_{j}=G_{j} H_{j}$, and let $S=\left[S_{0}, \ldots, S_{p}\right]$ and $w=\left[\mathbf{z}_{0}^{T}, \ldots, \mathbf{z}_{p}^{T}\right]^{T}$. Then

$$
\begin{equation*}
\left\|\mathcal{F}(\mathbf{y})-\sum_{j=0}^{p} S_{j j}\right\|_{\Omega}^{2}=\operatorname{tr}\left\{\left(E_{x x}-E_{x w} S^{T}-S E_{w x}+S E_{w w} S^{T}\right)\right\} \tag{3.7}
\end{equation*}
$$

where $E_{x w}=\left[E_{x z_{0}}, \ldots, E_{x z_{p}}\right]$ and by Theorem $1, E_{w w}=\operatorname{diag}\left[E_{z_{0} z_{0}}, \ldots, E_{z_{p} z_{p}}\right]$ is block-diagonal matrix. Thus,

$$
S E_{w w} S^{T}=S_{0} E_{z_{0} z_{0}} S_{0}^{T}+\ldots+S_{p} E_{z_{p} z_{p}} S_{p}^{T}
$$

$$
S E_{w x}=S_{0} E_{z_{0} x}+\ldots+S_{p} E_{z_{p} x}
$$

Therefore, (3.7) implies

$$
\begin{equation*}
\left\|\mathcal{F}(\mathbf{y})-\sum_{j=0}^{p} G_{j} H_{j j}\right\|_{\Omega}^{2}=\sum_{j=0}^{p}\left\|\mathcal{F}(\mathbf{y})-G_{j} H_{j} \mathbf{z}_{\mathbf{j}}\right\|_{\Omega}^{2}-\operatorname{tr}\left\{p E_{x x}\right\} \tag{3.8}
\end{equation*}
$$

Let us denote by $\mathbf{R}_{r}^{m \times n}$ the set of all $m \times n$ matrices of rank at most $r$. Then on the basis of $[83,30,52,12,87]$, the minimal Frobenius norm solution to the problem

$$
\begin{equation*}
\min _{G_{j} H_{j} \in \mathbf{R}_{r_{j}}^{m \times n}}\left\|\mathcal{F}(\mathbf{y})-G_{j} H_{j} \mathbf{z}_{j}\right\|_{\Omega}^{2} \tag{3.9}
\end{equation*}
$$

for $j=0, \ldots, p$, is given by

$$
\begin{equation*}
G_{j} H_{j}=U_{\Gamma_{z_{j}}, r_{j}} U_{\Gamma_{z_{j}}, r_{j}}^{T} E_{x z_{j}} E_{z_{j} z_{j}}^{\dagger} \tag{3.10}
\end{equation*}
$$

Therefore, (3.6) follows from (3.10).
Remark 1. Definition 3 of the well-defined injections is motivated by the following observation. It follows from (3.6) that if, for all $j=0, \ldots, p$, vector $\mathbf{v}_{j}$ is such that $\Gamma_{z_{j}}=\mathbf{O}$, then $G_{j}=\mathbf{O}$ and $H_{j}=\mathbf{O}$. In other words, then approximating operator $\mathcal{T}_{p}=\mathbf{O}$.

Therefore, in Theorem 2 above and in the Theorems below, vectors $\mathbf{v}_{0}, \ldots, \mathbf{v}_{p}$ are assumed well-defined.

### 3.3. Empirical determination of well-defined injections

First, we wish to show that the proposed system model $\mathcal{T}_{p}$ is also applicable under the conventional assumption used, for example, in $[15,36,65,83,35$, 70], when only matrices $E_{x x}, E_{x y}$ and $E_{y y}$ are assumed to be givenFootnoteIn particular, if $\mathbf{y}=A \mathbf{x}+\boldsymbol{\xi}$ where $A$ is a known matrix and random vector $\boldsymbol{\xi}$ is independent on $\mathbf{x}$, then it is assumed that only $E_{x x}$ and $E_{\mathbf{x} i, \mathbf{x} i}$ are given., and matrices $E_{x v_{j}}, E_{y v_{j}}, E_{v_{\ell} v_{j}}$ and associated joint probability functions are unknown. In this case, $E_{x v_{j}}, E_{y v_{j}}$ and $E_{v_{\ell} v_{j}}$ should be represented in terms of $E_{x x}, E_{x y}$ and $E_{y y}$.

We use the Stein's lemma [76] to show that it can be done indeed and, therefore, suppose that, for $q_{j} \geq n$, vector $\mathbf{y}=\left[\mathbf{y}_{(1)}, \ldots, \mathbf{y}_{(n)}\right]^{T}$ is extended to vector $\widetilde{\mathbf{y}}^{(j)}$ with $q_{j}$ components, $\widetilde{\mathbf{y}}^{(j)}=\left[\widetilde{\mathbf{y}}_{1}^{(j)}, \ldots, \widetilde{\mathbf{y}}_{q_{j}}^{(j)}\right]^{T}$, where $s_{k}$ entries of $\widetilde{\mathbf{y}}^{(j)}$ coincide with $\mathbf{y}_{(k)}$, for $k=1, \ldots, n$, and $q_{j}=$
$s_{1}+\ldots+s_{n}$ FootnoteThere are $q_{j}!/ s_{1}!\ldots s_{n}!$ such orderings.. In fact, $s_{1}, \ldots, s_{n}$ depend on $j$ but we omit subscript $j$ for the notation simplicity. (For example, if $\mathbf{y}=\left[\mathbf{y}_{(1)}, \mathbf{y}_{(2)}\right]^{T}$ then one of such extensions is $\widetilde{\mathbf{y}}^{(j)}=\left[\widetilde{\mathbf{y}}_{1}^{(j)}, \widetilde{\mathbf{y}}_{2}^{(j)}, \widetilde{\mathbf{y}}_{3}^{(j)}\right]^{T}=\left[\mathbf{y}_{1}^{(j)}, \mathbf{y}_{2}^{(j)}, \mathbf{y}_{1}^{(j)}\right]^{T}$ where $s_{1}=2$ and $s_{2}=1$.)

In general, injection $\mathbf{v}_{j}$, for $j=1, \ldots, p$, is defined by a nonlinear transformation $f_{j}$ so that $\mathbf{v}_{j}=f_{j}(\mathbf{y}) \in L^{2}\left(\Omega, \mathbf{R}^{q_{j}}\right)$ and $q_{j} \geq n$. Therefore, $\mathbf{v}_{j}=f_{j}(\mathbf{y})$ can be represented as $\mathbf{v}_{j}=\varphi_{j}\left[\phi_{j}(\mathbf{y})\right]=\varphi_{j}\left[\widetilde{\mathbf{y}}^{(j)}\right]$ where $\widetilde{\mathbf{y}}^{(j)}=$ $\phi_{j}(\mathbf{y}), \mathbf{v}_{j}=\left[\mathbf{v}_{1}^{(j)}, \ldots, \mathbf{v}_{q_{j}}^{(j)}\right]^{T}$ and $\mathbf{v}_{k}^{(j)}=\varphi_{j k}\left(\widetilde{\mathbf{y}}_{k}^{(j)}\right)$, for $k=1, \ldots, q_{j}$, i.e., $\left[\mathbf{v}_{1}^{(j)}, \ldots, \mathbf{v}_{q_{j}}^{(j)}\right]^{T}=\left[\varphi_{j k}\left(\widetilde{\mathbf{y}}_{1}^{(j)}\right), \ldots, \varphi_{j k}\left(\widetilde{\mathbf{y}}_{q_{j}}^{(j)}\right)\right]^{T}$.

By the conditions of the Stein's lemma [76], we also suppose that
(A1) $\varphi_{j}$ is a smooth function,
(A2) $E\left[\left|\varphi_{j}^{(n)}\left(\widetilde{\mathbf{y}}_{j}\right)\right|\right]<\infty$,
(A3) $E\left[\varphi_{j}^{\prime}\left(\widetilde{\mathbf{y}}_{j}\right)\right] \neq 0$,
(A4) sequence $\left\{a_{n}\right\}_{n=0}^{\infty}$ where $a_{n}=E\left[\varphi_{i}^{(n)}\left(\widetilde{\mathbf{x}}_{i}\right)\right] \times E\left[\varphi_{j}^{(n)}\left(\widetilde{\mathbf{x}}_{j}\right)\right]$ is monotone and bounded, for all $i, j=1, \ldots, q$, and
(A5) $\quad \mathbf{y}_{i}$ and $\varphi_{j}\left(\widetilde{\mathbf{y}}_{k}^{(j)}\right)$ are jointly normally distributed, for all $i, j=$ $1, \ldots, q$.

Let us write $E_{y \widetilde{y}(j)}=\left\{E_{y_{i} \widetilde{y}_{k}^{(j)}}\right\}_{i, k=1}^{n, q_{j}}$. Then by the Stein's lemma,

$$
E_{y_{i} v_{k}^{(j)}}=E_{y_{i} \varphi_{j k}\left(\widetilde{\mathbf{y}}_{k}^{(j)}\right)}=E_{y_{i}{\widetilde{y_{k}}}_{(j)}} E\left[\varphi_{j k}^{\prime}\left(\widetilde{\mathbf{y}}_{k}^{(j)}\right)\right] .
$$

Therefore,

$$
\begin{equation*}
E_{y v_{j}}=E_{y \widetilde{y}^{(j)}} D_{j}, \tag{3.11}
\end{equation*}
$$

where

$$
D_{j}=\operatorname{diag}\left(E\left[\varphi_{j 1}^{\prime}\left(\widetilde{\mathbf{y}}_{1}^{(j)}\right)\right], \ldots, E\left[\varphi_{j q_{j}}^{\prime}\left(\widetilde{\mathbf{y}}_{q_{j}}^{(j)}\right)\right]\right) \neq \mathbf{O}
$$

Similarly,

$$
\begin{equation*}
E_{x v_{j}}=E_{x \widetilde{y}\left({ }^{(j)}\right.} D_{j} . \tag{3.12}
\end{equation*}
$$

Thus, $E_{y v_{j}} \neq \mathbf{O}$ and $E_{x v_{j}} \neq \mathbf{O}$. Now, to obtain $E_{v_{\ell} v_{j}}$, we need to compute each entry $E_{v_{i}^{(\ell)} v_{k}^{(j)}}=E_{\varphi_{\ell i}\left(\widetilde{y}_{i}^{(\ell)}\right)} \varphi_{j k}\left(\widetilde{y}_{k}^{(j)}\right)$, for $i, k=1,2, \ldots, q_{j}$. By the results obtained in Section 3 of [59],

$$
\begin{equation*}
E_{\left.\varphi \ell i\left(\widetilde{y}_{i}^{(j)}\right) \varphi_{j k} \widetilde{y}_{k}^{(j)}\right)}=\sum_{n=1}^{\infty} \frac{\left(E_{\widetilde{y}_{i}^{(\ell)} \widetilde{y}_{k}^{(j)}}\right)^{n}}{n!} E\left[\varphi_{\ell i}^{(n)}\left(\widetilde{\mathbf{y}}_{i}^{(\ell)}\right)\right] E\left[\varphi_{j k}^{(n)}\left(\widetilde{\mathbf{y}}_{k}^{(j)}\right)\right] . \tag{3.13}
\end{equation*}
$$

On the basis of assumption (A4) and the Abel's test (see, e.g., Corollary 5.58 in [64]), the series in (3.13) converges.

In (3.11), (3.12) and (3.13), entries of matrices $E_{\widetilde{y}^{(j)}}, E_{\widetilde{y}^{(j)}}$ and $E_{\widetilde{y}_{i}^{(j)} \widetilde{y}_{k}^{(j)}}$ are the same as entries of matrices $E_{y y}, E_{x y}$ and $E_{y y}$, respectively. Therefore, $E_{y v_{j}}, E_{x v_{j}}$ and $E_{v_{\ell} v_{k}}$ are determined, in fact, in terms of $E_{x y}$ and $E_{y y}$, and expectations $E\left[\varphi_{\ell i}^{\prime}\left(\widetilde{\mathbf{y}}_{i}^{(\ell)}\right)\right]$ and $E\left[\varphi_{j k}^{(n)}\left(\widetilde{\mathbf{y}}_{k}^{(j)}\right)\right]$, for $i, k=1, \ldots, q_{j}$.

The following Example 1 illustrates the above derivations. Injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ defined in Example 1, are well-defined in the sense of condition (3.1) checked up numerically.

Example 1. Random vectors of type $\left[e^{\widetilde{\mathbf{y}}_{1}}, \ldots, e^{\widetilde{\mathbf{y}}_{q}}\right]^{T}$ have been exploited, in particular, in [43]. In this regard, let us consider $\mathbf{v}_{k}^{(j)}=\varphi_{j k}\left(\widetilde{\mathbf{y}}_{k}^{(j)}\right)=e^{\widetilde{\mathbf{y}}_{k}^{(j)}}$, for $j=1, . ., p$ and $k=1, \ldots, q_{j}$, and suppose the above conditions (A1)(A5) hold. Then $E\left[\varphi_{j k}^{\prime}\left(\widetilde{\mathbf{y}}_{k}^{(j)}\right]=E\left[\varphi_{j k}\left(\widetilde{\mathbf{y}}_{k}^{(j)}\right]=E\left[e^{\widetilde{\mathbf{y}}_{k}^{(j)}}\right]=\sqrt{e}\right.\right.$. Therefore, $D_{j}=\sqrt{e} I$ and then $E_{y v_{j}}=\sqrt{e} E_{y \widetilde{y}}{ }^{(j)}$ and $E_{x v_{j}}=\sqrt{e} E_{x \widetilde{y}(j)}$. On the other hand, $E\left[\varphi_{j i}^{(n)}\left(\widetilde{\mathbf{y}}_{i}^{(j)}\right)\right]=E\left[\varphi_{j i}\left(\widetilde{\mathbf{y}}_{i}^{(j)}\right)\right]=\sqrt{e}$. Thus, (3.13) implies

$$
E_{v_{i}^{(\ell)} v_{k}^{(j)}}=E_{\varphi_{\ell i}\left(\widetilde{y}_{i}^{(\ell)}\right) \varphi_{j k}\left(\widetilde{y}_{k}^{(j)}\right)}=e \sum_{n=1}^{\infty} \frac{\left(E_{\widetilde{y}_{i}^{(\ell)} \widetilde{y}_{k}^{(j)}}\right)^{n}}{n!}=e^{E_{\dddot{y}_{i}^{(\ell)} \widetilde{y}_{k}^{(j)}}+1}-e .
$$

Therefore, $E_{v_{\ell} v_{k}}=e M_{\ell k}$ where $M_{\ell k}=\left\{e^{E_{\tilde{y}_{i}(\ell)}^{y_{y}^{(j)}}}-1\right\}_{i, k=1}^{q_{j}}$. Then, in particular, for $j=1, E_{x z_{1}}=E_{x v_{1}}-E_{x y} E_{y y}^{\dagger} E_{y v_{1}}=\sqrt{e}\left(E_{x \widetilde{y}^{(1)}}-E_{x y} E_{y y}^{\dagger} E_{y \widetilde{y}(1)}\right)$ and $E_{z_{1} z_{1}}=e\left(M_{11}-E_{\widetilde{y}^{(1)} y} E_{y y}^{\dagger} E_{y \widetilde{y}^{(1)}}\right)$, where $\widetilde{\mathbf{y}}^{(1)} \neq \mathbf{y}$. That is, condition (3.1) is represented in terms of entries of matrices $E_{x y}$ and $E_{y y}$, as supposed, and can be verified numerically.

Remark 2. In our simulations in Examples 2, 3 and 4 below, empirical injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ are chosen 'similar' to vector $\mathbf{x}$. That is, if by the assumption, $\mathbf{x}$ is, e.g., Gaussian then $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ are Gaussian as well. This is another way to define $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$.

### 3.4. Error analysis associated with solution of problem (2.6)-(2.8)

Here, in Theorem 3, we obtain the constructive representation of the error associated with the solution of problem (2.6)-(2.8). We also prove that the error can be improved by the increase in the dimensions of injections
$\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ (Theorem 4) and degree of approximating operator $\mathcal{T}_{p}$ (Theorems $3,4,5)$.

Let us write

$$
\varepsilon_{G H}^{(p)}=\min _{G_{0}, H_{0}, \ldots, G_{p}, H_{p}}\left\|\mathcal{F}(\mathbf{y})-\sum_{j=0}^{p} G_{j} H_{j j}\right\|_{\Omega}^{2} .
$$

Theorem 3. For $j=0, \ldots, p$, let $A_{j}=E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}$, $\operatorname{rank}\left(A_{j}\right)=s_{j}$ and $s_{j} \geq r_{j}+1$. For $k=1, \ldots, s_{j}$, let $\sigma_{k}$ be a singular value of $A_{j}$. Then the error associated with the solution of problem (2.6)-(2.8) is represented by

$$
\begin{equation*}
\varepsilon_{G H}^{(p)}=\operatorname{tr}\left\{E_{x x}\right\}-\sum_{j=0}^{p} \sum_{k=1}^{r_{j}} \sigma_{k}^{2}\left(A_{j}\right) \tag{3.14}
\end{equation*}
$$

In particular, the error decreases as $p$ increases.
Proof. For $G_{j} H_{j}$ determined by (3.10), we have

$$
\begin{aligned}
\left\|\mathcal{F}(\mathbf{y})-G_{j} H_{j} \mathbf{z}_{\mathbf{j}}\right\|_{\Omega}^{2} & =\operatorname{tr}\left\{E_{x x}-E_{x z_{j}}\left(G_{j} H_{j}\right)^{T}-G_{j} H_{j} E_{z_{j} x}+G_{j} H_{j} E_{z_{j} z_{j}}\left(G_{j} H_{j}\right)^{T}\right\} \\
& =\left\|E_{x x}^{1 / 2}\right\|^{2}-\left\|E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2}+\left\|\left(G_{j} H_{j}-E_{x z_{j}} E_{z_{j} z_{j}}^{\dagger}\right) E_{z_{j} z_{j}}^{1 / 2}\right\|^{2} \\
(3.15) & =\left\|E_{x x}^{1 / 2}\right\|^{2}-\left\|E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2}+\left\|G_{j} H_{j} E_{z_{j} z_{j}}^{1 / 2}-E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2}
\end{aligned}
$$

because

$$
E_{z_{j} z_{j}}^{\dagger} E_{z_{j} z_{j}}^{1 / 2}=\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}
$$

and

$$
\begin{equation*}
E_{x z_{j}} E_{z_{j} z_{j}}^{\dagger} E_{z_{j} z_{j}}=E_{x z_{j}} \tag{3.16}
\end{equation*}
$$

(see [80]). Further, it follows from (3.6) that in the notation introduced in (3.5), matrix $G_{j} H_{j}$ is represented as $G_{j} H_{j}=\left[E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right]_{r_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right.$. Therefore,

$$
\begin{align*}
\left\|G_{j} H_{j} E_{z_{j} z_{j}}^{1 / 2}-E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2} & =\left\|\left[E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right]_{r_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger} E_{z_{j} z_{j}}^{1 / 2}-E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2} \\
& =\left\|\left[E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right]_{r_{j}}-E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2} \\
& =\left\|\left[A_{j}\right]_{r_{j}}-A_{j}\right\|^{2} \\
& =\sum_{k=r_{j}+1}^{s_{j}} \sigma_{k}^{2}\left(A_{j}\right) . \tag{3.17}
\end{align*}
$$

because $\left[E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right]_{r_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger} E_{z_{j} z_{j}}^{1 / 2}=\left[E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right]_{r_{j}}$. Further, since

$$
\begin{equation*}
\left\|E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2}=\left\|A_{j}\right\|^{2}=\sum_{k=1}^{s_{j}} \sigma_{k}^{2}\left(A_{j}\right) \tag{3.18}
\end{equation*}
$$

then (3.8), (3.15), (3.17) and (3.18) imply (3.14). In particular, it follows from (3.14) that the error decreases as degree $p$ increases.

Remark 3. According to Remark 2, in Examples 2, 3 and 4 that follow, empirical injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ are chosen 'similar' to vector $\mathbf{x}$. They are well-defined in the sense of condition (3.1) verified numerically.

Further, we wish to show the utility of injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ in the sense of the diminution of the associated error by increasing dimensions of $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$. To this end, we write

$$
A_{j}=\left\{a_{k i(j)}\right\}_{k, i=1}^{m, q_{j}} \text { and } A_{j}-\left[A_{j}\right]_{r_{j}}=\left\{b_{k i(j)}\right\}_{k, i=1}^{m, q_{j}}
$$

where $a_{k i(j)}$ and $a_{k i(j)}$ are entries of matrices $A_{j}$ and $A_{j}-\left[A_{j}\right]_{r_{j}}$, respectively. Let us also denote

$$
\begin{gathered}
\gamma_{k,(j)}=\sum_{i=1}^{m}\left(a_{k i(j)}^{2}-b_{k i(j)}^{2}\right), \\
\gamma_{(j)}=\max \left\{\gamma_{1,(j)}, \ldots, \gamma_{q_{j},(j)}\right\}, \\
\gamma=\max _{j=1, \ldots, p} \gamma_{(j)},
\end{gathered}
$$

$\alpha_{0}=\operatorname{tr}\left\{E_{x x}-\sum_{j=0}^{p} A_{j} A_{j}^{T}\right\}$, and $q=q_{1}+\ldots,+q_{p}$. We wish to show the utility of injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ in the sense of the diminution of the associated error by increasing dimensions of $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$.

Theorem 4. Let $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ be well-defined injections and let matrices $G_{0}, H_{0}, \ldots, G_{p}, H_{p}$ be defined by Theorem 2. Then the associated error decreases as the sum $q$ of dimensions of injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ increases. In particular, there is $\beta \in(0, \gamma]$ such that, given $\alpha \geq \alpha_{0}$, then

$$
\begin{equation*}
\alpha_{0} \leq \varepsilon_{G H}^{(p)} \leq \alpha \quad \text { iff } \quad q \geq \frac{\operatorname{tr}\left\{E_{x x}\right\}-\sum_{k=1}^{r_{0}} \sigma_{k}^{2}\left(A_{0}\right)}{\beta} . \tag{3.19}
\end{equation*}
$$

Proof. It follows from (3.8), (3.14), (3.15) and (3.17) that

$$
\begin{align*}
\varepsilon_{G H}^{(p)} & =\left\|\mathcal{F}(\mathbf{y})-G_{j} H_{0} \mathbf{z}_{\mathbf{0}}\right\|_{\Omega}^{2}+\sum_{j=1}^{p}\left\|\mathcal{F}(\mathbf{y})-G_{j} H_{j} \mathbf{z}_{\mathbf{j}}\right\|_{\Omega}^{2}-\operatorname{tr}\left\{p E_{x x}\right\} \\
20) & =\operatorname{tr}\left\{E_{x x}\right\}-\sum_{k=1}^{r_{0}} \sigma_{k}^{2}\left(A_{0}\right)+\sum_{j=1}^{p}\left\|\mathcal{F}(\mathbf{y})-G_{j} H_{j} \mathbf{z}_{\mathbf{j}}\right\|_{\Omega}^{2}-\operatorname{tr}\left\{p E_{x x}\right\} \tag{3.20}
\end{align*}
$$

where

$$
\begin{aligned}
\left\|\mathcal{F}(\mathbf{y})-G_{j} H_{j} \mathbf{z}_{\mathbf{j}}\right\|_{\Omega}^{2} & =\operatorname{tr}\left\{E_{x x}\right\}-\left(\left\|A_{j}\right\|^{2}-\left\|\left[A_{j}\right]_{r_{j}}-A_{j}\right\|^{2}\right) \\
& =\operatorname{tr}\left\{E_{x x}\right\}-\sum_{k=1}^{q_{j}} \sum_{i=1}^{m}\left(a_{k i(j)}^{2}-b_{k i(j)}^{2}\right) .
\end{aligned}
$$

Therefore,

$$
\begin{equation*}
\varepsilon_{G H}^{(p)}=\operatorname{tr}\left\{E_{x x}\right\}-\sum_{k=1}^{r_{0}} \sigma_{k}^{2}\left(A_{0}\right)-\sum_{j=1}^{p} \sum_{k=1}^{q_{j}} \sum_{i=1}^{m}\left(a_{k i(j)}^{2}-b_{k i(j)}^{2}\right) . \tag{3.21}
\end{equation*}
$$

Here, $\sum_{k=1}^{q_{j}} \sum_{i=1}^{m}\left(a_{k i(j)}^{2}-b_{k i(j)}^{2}\right)>0$ since by (3.17),

$$
\begin{equation*}
\left\|A_{j}\right\|^{2}-\left\|\left[A_{j}\right]_{r_{j}}-A_{j}\right\|^{2}=\sum_{k=1}^{r_{j}} \sigma_{k}^{2}\left(A_{j}\right)>0 \tag{3.22}
\end{equation*}
$$

Thus, (3.20) - (3.22) imply that $\varepsilon_{G H}^{(p)}$ decreases as $q_{j}$ increases, for $j=$ $1, \ldots, p$, and $p$ increases. Further, (3.14) implies

$$
\begin{equation*}
\varepsilon_{G H}^{(p)} \geq \operatorname{tr}\left\{E_{x x}\right\}-\sum_{j=0}^{p} \sum_{k=1}^{s_{j}} \sigma_{k}^{2}\left(A_{j}\right)=\operatorname{tr}\left\{E_{x x}-\sum_{j=0}^{p} A_{j} A_{j}^{T}\right\}=\alpha_{0} . \tag{3.23}
\end{equation*}
$$

Since

$$
0<\sum_{j=1}^{p} \sum_{k=1}^{q_{j}} \sum_{i=1}^{m}\left(a_{k i(j)}^{2}-b_{k i(j)}^{2}\right) \leq \gamma \sum_{j=1}^{p} q_{j}=\gamma q,
$$

then

$$
\begin{equation*}
\sum_{j=1}^{p} \sum_{k=1}^{q_{j}} \sum_{i=1}^{m}\left(a_{k i(j)}^{2}-b_{k i(j)}^{2}\right)=q \beta . \tag{3.24}
\end{equation*}
$$

Therefore, (3.21), (3.23) and (3.24) imply

$$
\alpha_{0} \leq \varepsilon_{G H}^{(p)}=\operatorname{tr}\left\{E_{x x}\right\}-\sum_{k=1}^{r_{0}} \sigma_{k}^{2}\left(A_{0}\right)-q \beta
$$

thus, if $\varepsilon_{G H}^{(p)} \leq \alpha$ then the right inequality in (3.19) is true. Conversely, if the latter is true then $\varepsilon_{G H}^{(p)} \leq \alpha$.

Remark 4. An empirical explanation of Theorem 4 is that the increase in $q$ implies the increase in the dimensions of matrices $H_{1}, \ldots, H_{p}$ in (2.5) and (3.6). Hence, it implies the increase in the number of parameters to optimize. As a result, for the given $r$, the accuracy associated with the system model $\mathcal{T}_{p}$ improves. Further, it follows from (3.19) that, as $q$ increases, $\varepsilon_{G H}^{(p)}$ tends to $\alpha_{0}$ which is the error associated with the full rank model $\mathcal{T}_{h}$ (see (3.32)). The error is given by (3.36).

Remark 5. By Theorem 3, the error associated with solution of problem (2.6) decreases as degree $p$ of the system model $\mathcal{T}_{p}$ increases. At the same time, the increase in degree $p$ of model $\mathcal{T}_{p}$ may involve an increase in parameter $r$ given by (2.4). However, by a condition of some applied problem in hand, $r$ must be fixed. In the following Theorem 5, under the condition of fixed $r$, the case of decreasing the error as the degree $p$ of the system model increases is detailed.

Theorem 5. Let $r$ and $r_{j}$, for $j=0, \ldots, p$, be given. Let $g$ be a nonnegative integer such that $g<p$ and let $\ell_{g}=r_{g}+r_{g+1}+\ldots+r_{p}$. If

$$
\begin{equation*}
\sum_{k=r_{g}+1}^{r_{g+1}+\ldots+r_{p}} \sigma_{k}^{2}\left(A_{g}\right)<\sum_{j=g+1}^{p} \sum_{k=1}^{r_{j}} \sigma_{k}^{2}\left(A_{j}\right), \tag{3.25}
\end{equation*}
$$

where

$$
\sum_{j=g+1}^{p} \sum_{k=1}^{r_{j}} \sigma_{k}^{2}\left(A_{j}\right)=\sum_{j=g+1}^{p} \sum_{k=1}^{q_{j}} \sum_{i=1}^{m}\left(a_{k i(j)}^{2}-b_{k i(j)}^{2}\right),
$$

then

$$
\begin{equation*}
\varepsilon_{G H}^{(p)}<\varepsilon_{G H}^{(g)}, \tag{3.26}
\end{equation*}
$$

i.e., for the same $r$, the error associated with the higher degree model $p$ is less than that associated with the lower degree model $g$.

Proof. We write $r=r_{0}+\ldots+r_{g-1}+\ell_{g}$. Then

$$
\begin{align*}
\varepsilon_{G H}^{(g)} & =\operatorname{tr}\left\{E_{x x}\right\}-\sum_{j=0}^{g-1} \sum_{k=1}^{r_{j}} \sigma_{k}^{2}\left(A_{j}\right)-\sum_{k=1}^{\ell_{g}} \sigma_{k}^{2}\left(A_{g}\right) \\
.27) & =\operatorname{tr}\left\{E_{x x}\right\}-\sum_{j=0}^{g-1} \sum_{k=1}^{r_{j}} \sigma_{k}^{2}\left(A_{j}\right)-\sum_{k=1}^{r_{g}} \sigma_{k}^{2}\left(A_{g}\right)-\sum_{k=r_{g}+1}^{r_{g+1}+\ldots+r_{p}} \sigma_{k}^{2}\left(A_{g}\right) . \tag{3.27}
\end{align*}
$$

Thus, (3.14) and (3.27) imply (3.25) and (3.26).
Remark 6. The RHS in (3.25) increases as the dimension $q_{j}$ of at least single injection $\mathbf{v}_{j}$, for $j=g+1, \ldots, p$, increases while the LHS does not depend on $q_{j}$. In other words, one can always find $q_{j}$, for $j=g+1, \ldots, p$, such that the inequality in (3.25) is true.

Remark 7. If $p=1$ and $m=q_{0}=q_{1}$, then proposed method is equivalent to the GBT2 [83]. For $m=q_{0}=q_{1}$ and $r_{0}=r_{1}=m / 4$, the computational load of the GBT2 is estimated as $57 \mathrm{~m}^{3}$. The computational load associated with computation of $G_{j}$ and $H_{j}$ by (3.6) is represented in Table 3.1.

Table 3.1: Estimate of computational load associated with computation of $G_{j}$ and $H_{j}$.

| Particular parts of (3.6) | Number of flops |
| :--- | :---: |
| 1. $E_{z_{j} z_{j}}^{\dagger}$ | $14 q_{j}^{3}$ |
| 2. $E_{x z_{j}} E_{j_{j} z_{j}}^{\dagger}$ | $2 m q_{j}^{2}$ |
| 3. $\Gamma_{z_{j}}=E_{x z_{j}} E_{z_{j} z_{j}}^{\dagger} E_{z_{j} x}$ | $2 m^{2} q_{j}$ |
| 4. $U_{\Gamma_{z_{j}}}$ from $S V D$ of $\Gamma_{z_{j}}$ | $12 m^{3}$ |
| 5. $H_{j}=U_{T_{z_{j}}, r_{j}}^{T} E_{x z_{j}} E_{z_{j} z_{j}}^{\dagger}$ | $2 m^{2} q_{j}+m^{2} r_{j}$ |

It follows from Table 3.1 that the total computational load of the proposed method is given by

$$
\sum_{j=0}^{p}\left(14 q_{j}^{3}+2 m q_{j}^{2}+4 m^{2} q_{j}+m^{2} r_{j}+12 m^{3}\right)
$$

In particular, if $m=q_{0}=\ldots=q_{p}$ and $r_{j}=m / 4$, for all $j=0,1, \ldots, p$, then the computational load of the proposed method is $32 \mathrm{pm}^{3}$. A lower error is provided by the proposed method at the expense of increased computational and/or storage costs what, we believe, is quite natural. In general, a better

Table 3.2: Numerical characterizations of models $\mathcal{T}_{0}\left(\mathbf{v}_{0}\right), \mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ and $\mathcal{T}_{2}\left(\mathbf{v}_{0}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)$ in Case 1

| Model | $q_{0}$ | $q_{1}$ | $q_{2}$ | $r_{0}$ | $r_{1}$ | $r_{2}$ | MSE |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\mathcal{T}_{0}\left(\mathbf{v}_{0}\right)$ | 100 | N/A | N/A | 50 | N/A | N/A | 8.30 |
| $\mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ | 100 | 25 | N/A | 25 | 25 | N/A | 7.93 |
| $\mathcal{T}_{2}\left(\mathbf{v}_{0}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)$ | 100 | 25 | 500 | 17 | 17 | 16 | 7.03 |

Table 3.3: Numerical characterizations of models $\mathcal{T}_{0}\left(\mathbf{v}_{0}\right), \mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ and $\mathcal{T}_{2}\left(\mathbf{v}_{0}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)$ in Case 2

| Model | $q_{0}$ | $q_{1}$ | $q_{2}$ | $r_{0}$ | $r_{1}$ | $r_{2}$ | MSE |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\mathcal{I}_{0}\left(\mathbf{v}_{0}\right)$ | 100 | N/A | N/A | 50 | N/A | N/A | 8.30 |
| $\mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ | 100 | 200 | N/A | 25 | 25 | N/A | 7.61 |
| $\mathcal{T}_{2}\left(\mathbf{v}_{0}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)$ | 100 | 200 | 500 | 17 | 17 | 16 | 6.28 |

result is normally achieved by some additional efforts. This observation is illustrated by the following example.

Example 2. Here, we wish to numerically illustrate Theorems 3, 4 and 5. To this end, we assume that $\mathbf{y}=\mathbf{x}+\boldsymbol{\xi}$ where $\mathbf{x} \in L^{2}\left(\Omega, R^{m}\right)$ and $\boldsymbol{\xi} \in L^{2}\left(\Omega, R^{m}\right)$ are uniformly and normally distributed random vectors, respectively. Injections $\mathbf{v}_{1} \in L^{2}\left(\Omega, R^{q_{1}}\right)$ and $\mathbf{v}_{2} \in L^{2}\left(\Omega, R^{q_{2}}\right)$ are here chosen as uniformly distributed random vectors. Random vector $\boldsymbol{\xi}$ simulates noise and is uncorrelated with $\mathbf{x}$ and $\mathbf{v}_{\mathbf{1}}, \mathbf{v}_{\mathbf{2}}$. Covariance matrices $E_{x v_{j}}, E_{v_{i} v_{j}}$, for $i, j=0,1,2$, are represented by $E_{x v_{j}}=\frac{1}{s} X V_{j}^{T}$ and $E_{v_{i} v_{j}}=\frac{1}{s} V_{i} V_{j}^{T}$ where $X \in \mathbf{R}^{m \times s}$ and $V_{j} \in \mathbf{R}^{q_{j} \times s}$ are sample matrices of $\mathbf{x}$ and $\mathbf{v}_{j}$, respectively, for $j=0,1,2$. We choose $m=100$ and $r=50$. Then models $\mathcal{T}_{0}\left(\mathbf{v}_{0}\right)$, $\mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ and $\mathcal{T}_{2}\left(\mathbf{v}_{0}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)$ are specified by choosing particular values of $q_{j}$ and $r_{j}$, for $j=0,1,2$, as shown in Tables 1 and 2.

In the Tables, Cases 1 and 2 for specific values of $q_{j}$ and $r_{j}$, for $j=0,1,2$, are considered. The values of the associated mean square errors (MSE) are also included. In Figure 3.1 the MSE are illustrated diagrammatically. It follows from Tables 1,2 and Figure 3.1 that, for the same $r$, the error associated with the proposed system model decreases if degree $p$ or the sum $q$ of the injection dimensions increases.


Figure 3.1: Example 2: Diagrams of the errors associated with $\mathcal{T}_{0}\left(\mathbf{v}_{0}\right)$, $\mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ and $\mathcal{T}_{2}\left(\mathbf{v}_{0}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)$.

Here, we wish to compare the error associated with the proposed technique (represented by $\varepsilon_{G H}^{(p)}$ in Theorem 3) with that of the method considered in [83] and abbreviated as GBT2. Recall that GBT2 is represented by $B(\mathbf{y}, g)=B_{1}\left(B_{2} \mathbf{y}+B_{3} g\right)$ where $B: L^{2}\left(\Omega, \mathbf{R}^{q_{0}}\right) \times L^{2}\left(\Omega, \mathbf{R}^{q_{1}}\right) \rightarrow L^{2}\left(\Omega, \mathbf{R}^{m}\right)$, $g \in L^{2}\left(\Omega, \mathbf{R}^{q_{1}}\right)$, and $B_{1}, B_{2}, B_{3}$ are linear operators defined by matrices $B_{1} \in R^{m \times m}$ and $B_{j} \in R^{m \times q_{j-2}}$, for $j=2,3$. GBT2 is not assigned for a system modeling and for the problem in (2.6), (2.7) and (2.8). At the same time, in terms of the model degree introduced in Section 2.2, GBT2 can be interpreted as a 'second' degree model. Therefore, it is interesting to compare the associated errors. Let $\mathbf{u}=\left[\mathbf{y}^{T} g^{T}\right]^{T} \in L^{2}\left(\Omega, \mathbf{R}^{2 n}\right)$ and denote by $\varepsilon_{G B T 2}$ the error associated with GBT2.

Theorem 6. Let $A_{u}=E_{x u}\left(E_{u u}^{1 / 2}\right)^{\dagger}$. If

$$
\begin{equation*}
\sum_{j=0}^{p} \sum_{k=1}^{r_{j}} \sigma_{k}^{2}\left(A_{j}\right)>\sum_{\ell=1}^{r} \sigma_{\ell}^{2}\left(A_{u}\right) \tag{3.28}
\end{equation*}
$$

then

$$
\begin{equation*}
\varepsilon_{G H}^{(p)}<\varepsilon_{G B T 2} . \tag{3.29}
\end{equation*}
$$

Proof. By [83], $\varepsilon_{G B T 2}=\operatorname{tr}\left\{E_{x x}\right\}-\sum_{\ell=1}^{r} \sigma_{\ell}^{2}\left(A_{u}\right)$. Then the latter and (3.14) involve the theorem validity.

In other words, under the condition in (3.28), the error associated with the proposed technique is less than that of GBT2. Note that the condition in (3.28) is easy to satisfy by the increase in $p$, for example.

Remark 8. The following Example 3 numerically illustrates the theoretical results obtained above and also demonstrates how the proposed technique is extended to the case of complex-valued random vectors.

Example 3. Let $\mathbf{x}_{c} \in L^{2}\left(\Omega, \mathbf{C}^{m}\right)$ and $\mathbf{y}_{c} \in L^{2}\left(\Omega, \mathbf{C}^{m}\right)$, i.e., let $\mathbf{x}_{c}$ and $\mathbf{y}_{c}$ be random vectors with complex values. Then $\mathbf{x}_{c}=\mathbf{x}_{r}+j \mathbf{x}_{i}$ and $\mathbf{y}_{c}=$ $\mathbf{y}_{r}+j \mathbf{y}_{i}$ where $j=\sqrt{-1}$, and $\mathbf{x}_{r}, \mathbf{y}_{r} \in L^{2}\left(\Omega, \mathbf{R}^{m}\right)$ and $\mathbf{x}_{i}, \mathbf{y}_{i} \in L^{2}\left(\Omega, \mathbf{R}^{m}\right)$ are the real and imaginary parts of $\mathbf{x}_{c}$ and $\mathbf{y}_{c}$, respectively, represented by Gaussian random vectors with zero-mean and standard deviation one (abbreviated as $G(0,1)$ ). Let us consider the case when $\mathbf{y}_{c}$ is represented by

$$
\begin{equation*}
\mathbf{y}_{c}=A_{c} \mathbf{x}_{c}+\boldsymbol{\xi}_{c}, \tag{3.30}
\end{equation*}
$$

where $A_{c}=\operatorname{diag}\left(e^{j \phi_{1}}, e^{j \phi_{2}}, \ldots, e^{j \phi_{m}}\right) \in \mathbf{C}^{m \times m}, j=\sqrt{-1}, \phi_{i}$ is a random number generated from Gaussian distribution with zero-mean and standard deviation one, $\boldsymbol{\xi}_{c} \in L^{2}\left(\Omega, \mathbf{C}^{m}\right)$ with $E_{\mathbf{x} i_{c} \mathbf{x} i_{c}}=\sigma^{2} I$ and $\sigma \in R$. The random vectors of this type are considered in a number of applications such as those in $[73,49,53,33,58]$.

As shown, in particular, in [61, 1], the problems concerning complex random vector $\mathbf{x}_{c}$ can always be treated in terms of the associated real random vector

$$
\mathbf{x}=\left[\begin{array}{c}
\mathbf{x}_{r} \\
\mathbf{x}_{i}
\end{array}\right] \in L^{2}\left(\Omega, \mathbf{R}^{2 m}\right) .
$$

Let $\boldsymbol{\xi}_{c}=\boldsymbol{\xi}_{r}+j \boldsymbol{\xi}_{i}$ and $A_{c}=A_{r}+j A_{i}$. The complex-valued vector given by (3.30) holds if and only if the real-valued random vector $\mathbf{y}=A \mathbf{x}+\boldsymbol{\xi}$ holds, where

$$
\mathbf{y}=\left[\begin{array}{c}
\mathbf{y}_{r} \\
\mathbf{y}_{i}
\end{array}\right] \in L^{2}\left(\Omega, \mathbf{R}^{2 m}\right), \quad \boldsymbol{\xi}=\left[\begin{array}{c}
\boldsymbol{\xi}_{r} \\
\boldsymbol{\xi}_{i}
\end{array}\right] \in L^{2}\left(\Omega, \mathbf{R}^{2 m}\right)
$$

and

$$
A=\left[\begin{array}{cc}
A_{r} & -A_{i} \\
A_{i} & A_{r}
\end{array}\right] \in \mathbf{R}^{2 m \times 2 m} .
$$

See, for instance, Appendix II in [61]. The above implies

$$
E_{x x}=\left[\begin{array}{ll}
E_{x_{r} x_{r}} & E_{x_{r} x_{i}} \\
E_{x_{i} x_{r}} & E_{x_{i} x_{i}}
\end{array}\right] \in \mathbf{R}^{2 m \times 2 m}
$$

and $E_{\mathbf{x} i \mathbf{x} i}=\sigma^{2} I$. Further, we denote $E_{x x}=\left\{e_{k, \ell}\right\}_{k, \ell=1}^{2 m}$. In applications such as those associated with synthetic aperture radar interferometry (see, e.g., [53, 33]) matrix $E_{x x}$ has a so-called triangular-time-Gaussian shape generated by

$$
e_{k, \ell}=\left\{\begin{array}{cl}
1-\frac{|k-\ell|}{m-1} b_{k}, & \text { for }|k-\ell| \leq \frac{m-1}{b_{k}}  \tag{3.31}\\
0, & \text { otherwise },
\end{array}\right.
$$

where $b_{k}$ is a parameter (called in [53, 33] the normalized baseline parameter). Here, $E_{x x}$ is considered in this form.



Figure 3.2: Example 3: Diagrams of the errors associated with $\mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$, $\mathcal{T}_{0}\left(\mathbf{v}_{0}\right)$ and GBT2 [83].

In this example, the injection $\mathbf{v}_{1}$ is represented as complex random vector $\mathbf{v}_{1 c}=\mathbf{v}_{1 r}+j \mathbf{v}_{1 i} \in L^{2}\left(\Omega, \mathbf{C}^{q_{1}}\right)$, where $\mathbf{v}_{1 r} \in L^{2}\left(\Omega, \mathbf{R}^{q_{1}}\right)$ and $\mathbf{v}_{1 i} \in$ $L^{2}\left(\Omega, \mathbf{R}^{q_{1}}\right)$ are the $G(0,1)$. It is assumed that $\mathbf{v}_{1 c}$ is uncorrelated with $\boldsymbol{\xi}_{c}$. Further, we write

$$
\mathbf{v}_{1}=\left[\begin{array}{c}
\mathbf{v}_{1 r} \\
\mathbf{v}_{1 i}
\end{array}\right] \in L^{2}\left(\Omega, \mathbf{R}^{2 q_{1}}\right), \quad \mathbf{u}=\left[\begin{array}{c}
\mathbf{x} \\
\mathbf{v}_{1}
\end{array}\right] \in L^{2}\left(\Omega, \mathbf{R}^{2\left(m+q_{1}\right)}\right),
$$

and choose $\mathbf{v}_{1}$ such that covariance matrix

$$
E_{u u}=\left[\begin{array}{cc}
E_{x x} & E_{x v_{1}} \\
E_{v_{1} x} & E_{v_{1} v_{1}}
\end{array}\right] \in \mathbf{R}^{2\left(m+q_{1}\right) \times 2\left(m+q_{1}\right)}
$$

has a triangular-time-Gaussian shape. As a result, matrices $E_{x v_{1}}$ and $E_{v_{1} v_{1}}$ are determined as blocks of $E_{u u}$. Then it follows that $E_{v_{1} y}=E_{v_{1} x} A^{T}$ and $E_{y y}=A E_{x x} A^{T}+\sigma^{2} I$. Then $G_{0}, H_{0}, \ldots, G_{p}, H_{p}$ are determined by (3.6). In this case, condition (3.1) is true, i.e. $\mathbf{v}_{1}$ is well-defined.

We wish to compare the accuracy associated with $\mathcal{T}_{p}\left(\mathbf{v}_{0}, \ldots, \mathbf{v}_{p}\right)$ where $G_{0}, H_{0}, \ldots, G_{p}, H_{p}$ are determined by (3.6), for $p=0$ and $p=1$. We also wish to compare it with the accuracy associated with the method considered in [83] and called GBT2.

In Figure 3.2, for $m=10$ and $b_{k}=0,4$ diagrams of the mean square error (MSE) associated with $\mathcal{T}_{0}\left(\mathbf{v}_{0}\right), \mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ and GBT2 are represented as follows. Figure 3.2 (a): the errors versus dimension $q_{1}$ of injection $\mathbf{v}_{1}$. Figure 3.2 (b): the errors versus parameter $\sigma$ in noise $\boldsymbol{\xi}$. Figure 3.2 (c): the errors versus the reduction ratio.

It follows from Figure 3.2 that the error associated with $\mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ is less than those of $\mathcal{T}_{0}\left(\mathbf{v}_{0}\right)$ and GBT2. In particular, in Figure 3.1 (a), the diagram of MSE associated with $\mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ illustrates Theorems 4 and 6. First, the MSE decreases as $q$ increases as Theorem 4 states. Second, for $q \geq 30$, condition (3.28) of Theorem 6 is true and the MSE associated with $\mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ is less than those of GBT2. For $q<30$, the MSE associated with $\mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ is slightly equal than that of GBT2. This is because, for $q<30$, condition (3.28) is not satisfied.

### 3.5. Optimal model of the system with no reduction of input dimensionality

An important case of problem (2.6)-(2.8) is when the matrix product $G_{j} H_{j}$, for $j=0, \ldots, p$, is replaced with a full rank matrix $P_{j} \in R^{m \times q_{j}}$. Then the input-output map $\mathcal{T}_{h}$ given by

$$
\begin{equation*}
\mathcal{T}_{h}\left(\mathbf{v}_{0}, \ldots \mathbf{v}_{p}\right)=\sum_{k=0}^{h} P_{k k} \tag{3.32}
\end{equation*}
$$

is called the full rank system model of degree $h$. In Theorems 7 and 8 below, the optimal full rank model of degree $h$ is determined and justified as a solution of problem (3.33). In problem (3.33) no constraint similar to that in (2.7) is imposed.

Theorem 7. Let vectors $\mathbf{z}_{0}, \ldots, \mathbf{z}_{p}$ be pairwise uncorrelated. Then the minimal Frobenius norm solution to the problem

$$
\begin{equation*}
\min _{P_{0}, \ldots, P_{h}}\left\|\mathcal{F}(\mathbf{y})-\sum_{k=0}^{h} P_{k k}\right\|_{\Omega}^{2} \tag{3.33}
\end{equation*}
$$

is given, for $k=0, \ldots, h$, by

$$
\begin{equation*}
P_{k}=E_{x z_{k}} E_{z_{k} z_{k}}^{\dagger} . \tag{3.34}
\end{equation*}
$$

Proof. Similar to (3.8),

$$
\begin{equation*}
\left\|\mathcal{F}(\mathbf{y})-\sum_{j=0}^{h} P_{j j}\right\|_{\Omega}^{2}=\sum_{j=0}^{h}\left\|\mathbf{x}-P_{j} \mathbf{z}_{\mathbf{j}}\right\|_{\Omega}^{2}-\operatorname{tr}\left\{h E_{x x}\right\} \tag{3.35}
\end{equation*}
$$

The minimal Frobenius norm solution to the problem

$$
\min _{P_{j}}\left\|\mathbf{x}-P_{j} \mathbf{z}_{\mathbf{j}}\right\|_{\Omega}^{2}
$$

is given by (3.34).
Theorem 8. Let $A_{j}=E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}$. The error associated with the minimal Frobenius norm solution to the problem in (3.33) is represented by

$$
\begin{equation*}
\min _{P_{0}, \ldots, P_{h}}\left\|\mathcal{F}(\mathbf{y})-\sum_{j=0}^{h} P_{j j}\right\|_{\Omega}^{2}=\operatorname{tr}\left\{E_{x x}\right\}-\sum_{j=0}^{h}\left\|A_{j}\right\|^{2} \tag{3.36}
\end{equation*}
$$

where $\|\cdot\|$ is the Frobenius norm.
Proof. For $P_{j}$ determined by (3.34),
$\left(\$ \mathcal{F}(\mathbf{y})-P_{j}\left\|_{\Omega}^{2}=\operatorname{tr}\left\{E_{x x}\right\}-\right\| E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\left\|^{2}+\right\| P_{j} E_{z_{j} z_{j}}^{1 / 2}-E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger} \|^{2}\right.$.
Here,

$$
\begin{aligned}
\left\|P_{j} E_{z_{j} z_{j}}^{1 / 2}-E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2} & =\left\|E_{x z_{j}} E_{z_{j} z_{j}}^{\dagger} E_{z_{j} z_{j}}^{1 / 2}-E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2} \\
& =\left\|E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}-E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2} \\
& =0 .
\end{aligned}
$$

That is,

$$
\begin{equation*}
\left\|\mathcal{F}(\mathbf{y})-P_{j j}\right\|_{\Omega}^{2}=\operatorname{tr}\left\{E_{x x}\right\}-\left\|E_{x z_{j}}\left(E_{z_{j} z_{j}}^{1 / 2}\right)^{\dagger}\right\|^{2} . \tag{3.38}
\end{equation*}
$$

Then (3.36) follows from (3.8) and (3.38).

Remark 9. It follows from Theorem 8 that the error associated with the system model determined by (3.32) and Theorem 7 decreases if degree $h$ increases.


Figure 3.3: Example 4: Diagrams of the errors associated with $\mathcal{T}_{0}\left(\mathbf{v}_{0}\right)$, $\mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ and $\mathcal{T}_{2}\left(\mathbf{v}_{0}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)$.

Example 4. Here, we wish to numerically illustrate Remark 9 and Theorem 8. To this end, we consider $\mathbf{y}=\mathbf{x}+\boldsymbol{\xi}$ where $\mathbf{x} \in L^{2}\left(\Omega, R^{m}\right)$ and $\boldsymbol{\xi} \in L^{2}\left(\Omega, R^{m}\right)$ are uniformly and normally distributed random vectors, respectively, with $m=50$. Injections $\mathbf{v}_{1} \in L^{2}\left(\Omega, R^{q_{1}}\right)$ and $\mathbf{v}_{2} \in L^{2}\left(\Omega, R^{q_{2}}\right)$ are here chosen as uniformly distributed random vectors with $q_{1}=q_{2}=$ $m=50$. Noise $\boldsymbol{\xi}$ is uncorrelated with $\mathbf{x}$ and $\mathbf{v}_{\mathbf{1}}, \mathbf{v}_{\mathbf{2}}$. The mean square errors (MSE) associated with $\mathcal{T}_{0}\left(\mathbf{v}_{0}\right), \mathcal{T}_{1}\left(\mathbf{v}_{0}, \mathbf{v}_{1}\right)$ and $\mathcal{T}_{2}\left(\mathbf{v}_{0}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)$ are 11.81, 9.52 and 7.95, respectively. They are diagrammatically shown in Figure 3.3, and illustrate Remark 9 and Theorem 8, i.e., the error associated with model $\mathcal{T}_{h}$ decreases if degree $h$ increases.

## 4. Conclusion and open problem

We have proposed and justified a new approach to the optimal modeling of non-linear system $\mathcal{F}$. It is assumed that its random input $\mathbf{y}$ and random output $\mathbf{x}$ are available and known in terms covariance matrices, and the input-output map $\mathcal{F}$ is unknown. The basic idea is to build the system model $\mathcal{I}_{p}$ as a sum of $p+1$ particular parts. The constructive determination of each part of the model is provided. The proposed model is optimal in the sense of minimization of the associated error. A significant ingredient of the proposed model is a set of so called injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ which allows us to further improve the accuracy of the system representation. The proposed model has several degrees of freedom to minimize the associated error.

An open associated problem is as follows: Given $\mathbf{x}, \mathbf{y}$ and $r_{0}, \ldots, r_{p}$, find matrices $G_{0}, H_{0}, \ldots, G_{p}, H_{p}$, injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ and transformations $Q_{0}, \ldots, Q_{p}$ that solve

$$
\begin{equation*}
\min _{\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}} \min _{G_{0}, H_{0}, \ldots, G_{p}, H_{p}}\left\|\mathcal{F}(\mathbf{y})-\sum_{j=0}^{p} G_{j} H_{j} \mathbf{z}_{j}\right\|_{\Omega}^{2} \tag{4.1}
\end{equation*}
$$

subject to

$$
\begin{equation*}
G_{j} \in R^{m \times r_{j}} \text { and } H_{j} \in R^{r_{j} \times q_{j}} \tag{4.2}
\end{equation*}
$$

$$
\begin{equation*}
E_{z_{i} z_{j}}=\mathbf{O}, \quad \text { for } i \neq j \tag{and}
\end{equation*}
$$

where $i, j=0, \ldots, p$.
An important difference from the problem in (2.6), (2.7) and (2.8) is that the problem in $(4.1),(4.2)$ and (4.3) has $p$ more unknowns to determine, they are $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$. It makes this problem quite difficult. We are planning to tackle to it in the near future. A possible solution device is based on the special iterative procedure such that, for each iteration loop, the best approximation problem is solved providing searched optimal matrices $G_{0}, H_{0}, \ldots, G_{p}, H_{p}$ and optimal injections $\mathbf{v}_{1}, \ldots, \mathbf{v}_{p}$ with the smallest associated error.

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